On solutions of the Schlesinger Equations in Terms of Θ -Functions

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Abstract

In this paper we construct explicit solutions and calculate the corresponding τ -function to the system of Schlesinger equations describing isomonodromy deformations of 2×2 matrix linear ordinary differential equation whose coefficients are rational functions with poles of the first order; in particular, in the case when the coefficients have four poles of the first order and the corresponding Schlesinger system reduces to the sixth Painlevé equation with the parameters 1/8, -1/8, 1/8, 3/8, our construction leads to a new representation of the general solution to this Painlevé equation obtained earlier by K. Okamoto and N. Hitchin, in terms of elliptic thetafunctions.

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1 Introduction

The Schlesinger equations [17] arise in the context of the following Riemann-Hilbert (inverse monodromy) problem:

for an arbitrary $g \in \mathbb{N}$ and distinct 2g + 2 points $\lambda_j \in \mathbb{C}$, construct a function $\Psi(\lambda)$: $\mathbb{CP}^1 \setminus \{\lambda_1, \dots, \lambda_{2g+2}\} \to \mathrm{SL}(2, \mathbb{C})$ which has the following properties;

- 1) $\Psi(\infty) = I$,
- 2) $\Psi(\lambda)$ is holomorphic for all $\lambda \in \mathbb{CP}^1 \setminus \{\lambda_1, \dots, \lambda_{2q+2}\}$,
- 3) $\Psi(\lambda)$ has regular singular points at $\lambda = \lambda_j$, $j = 1, \ldots, 2g + 2$, with given monodromy matrices, $M_j \in SL(2, \mathbb{C})$

In the case when the monodromy matrices are independent of the parameters $\lambda_1, \ldots, \lambda_{2g+2}$, the function $\Psi \equiv \Psi(\lambda)$ solves the following matrix differential equation,

$$\frac{d\Psi}{d\lambda} = \sum_{j=1}^{2g+2} \frac{A_j}{\lambda - \lambda_j} \Psi, \tag{1.1}$$

where the $sl(2,\mathbb{C})$ -valued matrices A_i solve the system of Schlesinger equations,

$$\frac{\partial A_j}{\partial \lambda_i} = \frac{[A_i, A_j]}{\lambda_i - \lambda_j}, \quad i \neq j , \qquad \frac{\partial A_i}{\partial \lambda_i} = -\sum_{j \neq i} \frac{[A_i, A_j]}{\lambda_i - \lambda_j}. \tag{1.2}$$

Obviously, the eigenvalues of A_j , which will be denoted by $\frac{t_j}{2}$ and $-\frac{t_j}{2}$ in the sequel, are integrals of motion of system (1.2).

The important object associated with system (1.2) is the so-called τ -function - the function generating Hamiltonians of the Schlesinger system [16, 8, 7]; it can by defined as the solution to the following system of equations,

$$\frac{\partial \ln \tau}{\partial \lambda_j} \equiv \sum_{i \neq j} \frac{\mathrm{tr} A_j A_i}{\lambda_j - \lambda_i}$$

(see Sec.2 for details).

For g=1, the Schlesinger system may equivalently be rewritten in terms of a single function of one variable, the position y(t) of the zero of the (12) matrix element of the function $\frac{A_1}{\lambda} + \frac{A_2}{\lambda - 1} + \frac{A_3}{\lambda - t}$ in the λ -plane. The equation for y(t) turns out to coincide with the sixth Painlevé equation,

$$\frac{d^2y}{dt^2} = \frac{1}{2} \left(\frac{1}{y} + \frac{1}{y-1} + \frac{1}{y-t} \right) \left(\frac{dy}{dt} \right)^2 - \left(\frac{1}{t} + \frac{1}{t-1} + \frac{1}{y-t} \right) \frac{dy}{dt} + \frac{y(y-1)(y-t)}{t^2(t-1)^2} \left(\alpha + \beta \frac{t}{y^2} + \gamma \frac{t-1}{(y-1)^2} + \delta \frac{t(t-1)}{(y-t)^2} \right),$$
(1.3)

where

$$\alpha \equiv \frac{(t_1 - 1)^2}{2}, \quad \beta \equiv -\frac{t_2^2}{2}, \quad \gamma \equiv \frac{t_3^2}{2}, \quad \delta \equiv \frac{1}{2} - \frac{t_4^2}{2},$$
 (1.4)

K. Okamoto showed [15] that the general solution to the sixth Painlevé equation can be written explicitly in terms of elliptic functions provided that the set of the parameters t_j satisfy one of the following conditions: $t_i \in \mathbb{Z}$, $t_1 + \ldots + t_4 \in 2\mathbb{Z}$ or $t_i + \frac{1}{2} \in \mathbb{Z}$. More recently, the algebro-geometric aspects of the sixth Painlevé equation have once again attracted the attention, see the papers [5, 13] (some details which are relevant to our work are given in Appendix).

Our interest to the problem of finding explicit solutions of the Schlesinger system in algebro geometric terms was initiated, on one hand, by the work of Okamoto, and, on the other hand, by our papers [10, 12, 11, 9], devoted to the study of solutions to the Ernst equation arising as a partial case of the vacuum Einstein equations; in particular, it turns out that some of the elliptic solutions of the Ernst equation studied in [11] may also be described by the sixth Painlevé equation [9]: in fact, being rewritten in appropriate variables, these solutions give a certain one-parameter sub-family of the Okamoto's solutions with $t_i = 1/2$.

In this paper we solve, in terms of theta-functions, the inverse monodromy problem formulated at the beginning of the Introduction for an arbitrary g and an arbitrary set of anti-diagonal monodromy matrices. Our approach originated from the so-called finite-gap integration method for the integrable systems [3]. The solution of the inverse monodromy problem allows, in turn, to express in terms of theta functions the 2g-parameter family of solutions to the Schlesinger system for $t_j = \frac{1}{2}$ and calculate the corresponding τ -function. In contrast to the common belief (which got its origin in the papers [7, 16, 8]) that for algebro geometrical solutions to integrable systems the τ -function simply coincides with certain theta functions, in the present case, the τ -function (up to multiplication by an arbitrary constant) is given by the following expression,

$$\tau(\{\lambda_j\}) = \frac{\Theta[\mathbf{p}, \mathbf{q}](0|\mathbf{B})}{\sqrt{\det \mathcal{A}}} \prod_{j \le k} (\lambda_j - \lambda_k)^{-\frac{1}{8}}, \tag{1.5}$$

where the vectors $\mathbf{p} \in \mathbb{C}^g$, $\mathbf{q} \in \mathbb{C}^g$ are parameters corresponding to parameters of the monodromy matrices, \mathbf{B} is the matrix of *b*-periods of the hyperelliptic curve

$$w^2 = \prod_{j=1}^{2g+2} (\lambda - \lambda_j),$$

and

$$\mathcal{A}_{kj} \equiv 2 \int_{\lambda_{2j+1}}^{\lambda_{2j+2}} \frac{\lambda^{k-1} d\lambda}{w} , \quad j, k = 1, \dots, g.$$

For the elliptic case g = 1, applying a conformal transformation of the λ -plane, one can always map the points $\lambda_1, \ldots, \lambda_4$ to 0, 1, t and ∞ , respectively (t is equal to the cross-ratio of the points $\lambda_1, \ldots, \lambda_4$). Then (again up to an arbitrary constant) the τ -function (1.5) can be rewritten in the following form,

$$\tau(t) = \frac{\theta_{p,q}(0|\sigma)}{\sqrt[8]{t(t-1)}} \left[\int_0^1 \frac{d\lambda}{\sqrt{\lambda(\lambda-1)(\lambda-t)}} \right]^{-\frac{1}{2}}, \tag{1.6}$$

where $\theta_{p,q}(0|\sigma)$ is the elliptic theta-function with characteristic [p,q]: here, the module $\sigma(t)$ of the curve $w^2 = \lambda(\lambda-1)(\lambda-t)$ is chosen so that $t = \theta_4^4(0|\sigma)/\theta_2^4(0|\sigma)$.

The latter τ -function defines a new representation of the solution to the sixth Painlevé equation with the parameters $t_i = 1/2$ i.e.

$$\alpha = \frac{1}{8}, \quad \beta = -\frac{1}{8}, \quad \gamma = \frac{1}{8}, \quad \delta = \frac{3}{8} :$$
 (1.7)

$$y(t) = t - t(t - 1) \left[D\left(\frac{\frac{d}{dt}D(\tau)}{\frac{d}{dt}D\left(\sqrt[8]{t(t - 1)}\tau\right)}\right) + \frac{t(t - 1)}{D^2\left(\sqrt[8]{t(t - 1)}\tau\right)} \right]^{-1}, \tag{1.8}$$

where the operator D is defined as follows,

$$D(\cdot) \equiv t(t-1)\frac{d}{dt}\ln(\cdot).$$

As a corollary of sixth Painlevé equation (1.3) with coefficients (1.7), function

$$\zeta(t) \equiv D(\tau)$$

where the τ -function $\tau(t)$ is given by (1.6), satisfies the following equation:

$$[t(t-1)\zeta'']^2 = \zeta'[(\zeta' + \frac{1}{4})^2 - ((2t-1)\zeta' - \zeta)^2]$$
(1.9)

One more form of the solution (1.8), namely,

$$y(t) = \frac{tu(\frac{\sigma}{2}|\sigma)}{u(\frac{\sigma}{2},|\sigma) + (1-t)u(\frac{1}{2}|\sigma)}, \quad \text{where} \quad u(z|\sigma) = \frac{\frac{\partial}{\partial z} \ln \frac{\partial}{\partial z} \ln \frac{\theta_{p,q}(z|\sigma)}{\theta_1(z|\sigma)}}{\frac{\partial}{\partial \sigma} \ln \frac{\theta_{p,q}(z|\sigma)}{\theta_1(z|\sigma)}}, \quad (1.10)$$

may be obtained from our construction by a straightforward calculation of the position of the zero of the (12) component of the matrix $\Psi_{\lambda}\Psi^{-1}$ in the λ -plane.

This paper is organized as follows. In Section 2, we recall some basic facts about isomonodromy deformations and Schlesinger equations. In Section 3, we begin with the solution of an inverse monodromy problem with an arbitrary even number of singular points and anti-diagonal monodromy matrices. In Section 4, we find the related τ -function, and finally, in Section 5, we apply the results of the previous sections to the g=1 case, i.e., to the sixth Painlevé equation.

It is also worth mentioning, that the solution of some inverse monodromy problems, including singularities of regular and irregular type in the framework of the finite-gap integration technique, were given by M. Jimbo and T. Miwa [7], however, their construction can not be applied to solve the inverse monodromy problems considered here. In the case of 2×2 monodromy problems with only regular singularities, say, the construction by Jimbo and Miwa leads to an analytic function with 3g + 2 regular singular points whose 2g + 2 monodromy matrices, after a proper normalization (see Section 2), equal $i\sigma_1$, and g monodromy matrices are just equal to -I. Therefore, the solution of the Schlesinger system, which can be obtained from the construction of Jimbo and Miwa, does not contain any parameters in contrast to the construction presented in this paper.

Simultaneously with the present work, solution of the same Riemann-Hilbert problem was given in the paper of P.Deift, A.Its, A.Kapaev and X.Zhou [2] in rather different terms. The problem of calculation of corresponding τ -function (1.5) was not considered there.

2 The Schlesinger Equations

In this section we recall the basic notations and definitions related to isomonodromy deformations of 2×2 matrix linear ordinary differential equation,

$$\frac{d}{d\lambda}\Psi = A(\lambda)\Psi,\tag{2.1}$$

where $A(\lambda) \in sl_2(\mathbb{C})$ is a rational function of λ with 2g + 2 poles of the first order,

$$A(\lambda) = \sum_{j=1}^{2g+2} \frac{A_j}{\lambda - \lambda_j}, \quad i \neq j \Rightarrow \lambda_i \neq \lambda_j, \quad \frac{d}{d\lambda} A_j = 0.$$
 (2.2)

We suppose that $\lambda = \infty$ is not a pole, which means that the following condition is fulfilled

$$\sum_{j=1}^{2g+2} A_j = 0. (2.3)$$

To fix a fundamental solution of Eq. (2.1), choose a point $\lambda_0 \in \mathbb{P} \setminus \{\lambda_1, \dots, \lambda_{2g+2}\}$ and impose the following normalization condition:

$$\Psi(\lambda_0) = I. \tag{2.4}$$

Since $\operatorname{tr} A(\lambda) = 0$, this means that $\det \Psi(\lambda) = 1$ for $\lambda \in \mathbb{C}$. Now one defines the monodromy matrices,

$$M_j = \Psi(\lambda_0)|_{\gamma_j}, \quad k = j, \dots, 2g + 2,$$

as analytic continuations of the fundamental solution normalized by condition (2.4) along the generators, γ_k , of the fundamental group $\pi_1(\mathbb{CP}^1 \setminus \{\lambda_1, \lambda_2, \dots, \lambda_{2g+2}\}, \lambda_0)$ defined in the figure 1.

The monodromy matrices satisfy the cyclic relation,

$$M_{2g+2} \cdot \ldots \cdot M_1 = I, \tag{2.5}$$

and generate a subgroup of $SL(2,\mathbb{C})$, i.e.,

$$\det M_i = 1, \quad i = 1, \dots, 2q + 2. \tag{2.6}$$

Matrix elements of M_j and eigenvalues $\pm \frac{t_j}{2}$ of the matrices A_j , $j=1,\ldots,2g+2$, are called the *monodromy data* of the function Ψ . The monodromy data are locally analytic functions of the variables A_1,\ldots,A_{2g+2} and $\lambda_0,\lambda_1,\ldots,\lambda_{2g+2}$. The condition

$$\frac{dt_j}{d\lambda_l} = 0$$
 and $\frac{dM_j}{d\lambda_l} = 0$, for $j, l = 1, \dots, 2g + 2$, (2.7)

is called the *isomonodromy condition*. The isomonodromy condition (2.7) is equivalent to the following system of linear differential equations for the function Ψ :

$$\frac{d\Psi}{d\lambda_j} = \left(\frac{A_j}{\lambda_0 - \lambda_j} - \frac{A_j}{\lambda - \lambda_j}\right)\Psi, \qquad j = 1, \dots, 2g + 2. \tag{2.8}$$

Following [17] we choose the normalization point $\lambda_0 = \infty$ to exclude the nonessential parameter λ_0 . In this case the compatibility condition of system (2.8), (2.1) reads as the following system of nonlinear ODEs, the *Schlesinger equations*:

$$j \neq i: \qquad \frac{\partial A_j}{\partial \lambda_i} = \frac{[A_i, A_j]}{\lambda_i - \lambda_j},$$
 (2.9)

$$j = i: \qquad \frac{\partial A_i}{\partial \lambda_i} = -\sum_{\substack{j=1\\i \neq i}}^{2g+2} \frac{[A_i, A_j]}{\lambda_i - \lambda_j}, \tag{2.10}$$

Solutions of these equations define isomonodromy deformations of the matrix elements of A_i . Note that system (2.9), (2.10) is equivalent to system (2.9), (2.3).

Proposition 2.1 If a set $\{A_1, \ldots, A_{2g+2}\}$ is a solution of the system (2.9), (2.10), then the monodromy data of the function Ψ , which solves Eq. (2.1) with the corresponding matrix $A(\lambda)$ given by Eq. (2.2), are independent of $\lambda_1, \ldots, \lambda_{2g+2}$.

The set of the monodromy data, $\{t_1, \ldots, t_{2g+2}, M_1, \ldots, M_{2g+2}\} \in \mathbb{C}^{2g+2} \times \mathcal{M}_{2g+2}$, where the variety $\mathcal{M}_{2g+2} \equiv \mathcal{M}_{2g+2}(t_1, \ldots, t_{2g+2})$ is defined via Eqs. (2.5) and (2.6), is known to be in one-to-one correspondence with the solutions of the system of Schlesinger equations (2.9), (2.10). The nontrivial part of this statement follows from the solvability of the inverse monodromy problem (see [1]).

In this paper we consider the case when all $t_j = 1/2$, so that the matrices A_j and M_j can be represented in the following form,

$$A_j = \frac{1}{4}G_j\sigma_3G_j^{-1}, \quad M_j = iC_j^{-1}\sigma_3C_j, \tag{2.11}$$

and λ -independent matrices G_j and C_j are defined via the asymptotic behavior of the function Ψ in the neighborhood of the points λ_j ,

$$\Psi = \underset{\lambda \to \lambda_j}{=} (G_j + \mathcal{O}(\lambda - \lambda_j))(\lambda - \lambda_j)^{\frac{1}{4}\sigma_3} C_j; \tag{2.12}$$

$$\det G_j = \det C_j = 1.$$

In the isomonodromy case one can always choose C_j to be independent of $\lambda_1, \ldots, \lambda_{2g+2}$. Hereafter we use the standard notation for the Pauli matrices:

$$\sigma_1 = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \quad \sigma_2 = \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix}, \quad \sigma_3 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}.$$

One can formulate the following

Proposition 2.2 Let $\Psi^*(Q)$ be a holomorphic function on the universal covering, $\operatorname{pr}: X \to \mathbb{CP}^1 \setminus \{\lambda_1, \dots, \lambda_{2g+2}\}$, which has the asymptotic behavior as $\lambda = \operatorname{pr} Q \to \lambda_j$ prescribed by Eq. (2.12) and normalized as $\Psi^*(Q_0) = I$ at some point Q_0 , $\operatorname{pr} Q_0 = \lambda_0$. Then the function $\Psi(\lambda) = \Psi^*(Q)|_{\operatorname{pr} Q = \lambda}$ has the monodromy data corresponding to the variety $\mathcal{M}_{2g+2}(\pm \frac{1}{2}, \dots, \pm \frac{1}{2})$, with the matrices M_j defined via the second equation (2.11), and solves the system of differential equations (2.1), (2.8), where the matrix $A(\lambda)$ is defined by Eqs. (2.1) and (2.2).

If a set of matrices $\{A_1, \ldots, A_{2g+2}\}$ is a solution of the system (2.9), (2.10), then for any matrix $K \in SL(2, \mathbb{C})$ independent of $\lambda_1, \ldots, \lambda_{2g+2}$ the new set $\{A_j^{new} = KA_jK^{-1}, j = 1, \ldots, 2g+2\}$ is also a solution of the system. This gauge transformation on the set of the solutions of the Schlesinger system corresponds to the following gauge transformation of the function $\Psi(\lambda)$,

$$\Psi^{new} = K\Psi K^{-1}.\tag{2.13}$$

which leaves the normalization condition (2.4) invariant and acts on \mathcal{M}_{2g+2} in the same way as on the space of the solutions,

$$M_j^{new} = K M_j K^{-1}.$$
 (2.14)

By choosing $K = C_0C_1$, where C_1 is given by (2.12) for j = 1 and $C_0 = \frac{i}{\sqrt{2}}(\sigma_3 + \sigma_1)$, we use this gauge transformation to fix

$$M_1 = i\sigma_1. (2.15)$$

Since we have one more parameter in our gauge transform, $C_0 \to C_0 \kappa^{\sigma_3}$, we can use the remaining freedom to remove one more parameter from \mathcal{M}_{2g+2} . More exactly, by making one more gauge transform (2.13) with the matrix $K = C_0 \kappa^{\sigma_3} C_0^{-1}$, we, by choosing appropriately the parameter κ , fix the next monodromy matrix M_2 :
If $\operatorname{tr}(M_2\sigma_1)^2 \neq -2$, then

$$M_2 = \begin{pmatrix} 0 & m_2 \\ -m_2^{-1} & 0 \end{pmatrix}, \quad m_2 \in \mathbb{C}^* \equiv \mathbb{C} \setminus \{0, \infty\}; \tag{2.16}$$

if $\operatorname{tr}(M_2\sigma_1)^2=-2$ but $M_2\neq \pm i\sigma_1$, then $M_2=\pm i(\sigma_3+\sigma_1+i\sigma_2)$; and, finally, if $M_2=\pm i\sigma_1$, then we can use the parameter κ to fix analogously the structure of the next matrix, M_3 .

The variety $\mathcal{M}_{2g+2}(\pm \frac{1}{2}, \dots, \pm \frac{1}{2})$ contains the following sub-variety, $\mathbb{C}_{2g}^* \cong \mathbb{T}_{2g} \times \mathbb{R}^{2g}$:

$$M_j = \begin{pmatrix} 0 & m_j \\ -m_j^{-1} & 0 \end{pmatrix}, \quad j = 1, \dots, 2g + 2,$$
 (2.17)

where

$$m_1 = i, \quad m_j \in \mathbb{C}^*, \quad j = 2, \dots, 2g + 2; \qquad \prod_{j=1}^{g+1} m_{2j} = (-1)^{g+1} \prod_{j=1}^{g+1} m_{2j-1}.$$
 (2.18)

Note that if the matrices M_1 and M_2 are fixed by Eqs. (2.15) and (2.16) correspondingly, then $\dim_{\mathbb{C}} \mathcal{M}_{2g+2}(\pm \frac{1}{2}, \dots, \pm \frac{1}{2}) = 4g - 2$ and $\dim_{\mathbb{C}} \mathbb{C}_{2g}^* = 2g$; in fact, for g = 1 the sub-variety \mathbb{C}_{2g}^* constitutes almost all the variety $\mathcal{M}_{2g+2}(\pm \frac{1}{2}, \dots, \pm \frac{1}{2})$. More precisely, one can formulate the following

Proposition 2.3 [5] If g = 1, then the variety $\mathcal{M}_4(\pm \frac{1}{2}, \dots, \pm \frac{1}{2})$ coincides, up to the conjugation defined by Eq. (2.14) with arbitrary matrix $K \in \mathrm{SL}(2,\mathbb{C})$, with the union of the following two sets of the monodromy matrices:

1)
$$M_k = \begin{pmatrix} 0 & m_k \\ -\frac{1}{m_k} & 0 \end{pmatrix}$$
, $k = 1, \dots 4$, $m_1 = i$, $m_k \in \mathbb{C}^*$, $m_4 m_2 = i m_3$; (2.19)

2)
$$M_1 = -i\sigma_3$$
, $M_2 = i\epsilon_2 \begin{pmatrix} -1 & a-1 \\ 0 & 1 \end{pmatrix}$, $M_3 = i\epsilon_3 \begin{pmatrix} -1 & a \\ 0 & 1 \end{pmatrix}$, $M_4 = i\epsilon_4 \begin{pmatrix} -1 & 1 \\ 0 & 1 \end{pmatrix}$, (2.20)

where $\epsilon_k = \pm 1$, $\epsilon_2 \epsilon_3 \epsilon_4 = 1$, $a \in \mathbb{C}$.

Isomonodromy deformations of Eq. (2.1) in the case when the matrix $A(\lambda)$ has four poles are governed by solutions to the sixth Painlevé equation (1.3). Here we rewrite the corresponding relation given by M. Jimbo and T. Miwa [7] in the notation which more suits to our basic construction:

Denote by \vec{g}_j^p the pth column of the matrix G_j from Eq.(2.11), and introduce new matrices $G_{ij}^{pq} \stackrel{\text{def}}{=} (\vec{g}_i^p \vec{g}_j^q)$; in particular, $G_{jj}^{12} \equiv G_j$.

Proposition 2.4 The functions

$$\hat{A}_{j}^{12} = t_{j} \frac{\det G_{j1}^{12} \det G_{1j}^{22}}{\det G_{11}^{12} \det G_{jj}^{12}}, \quad j = 1, \dots, 4,$$
(2.21)

depend on the variables $\{\lambda_k\}$ only through their cross-ratio,

$$t = \frac{\lambda_3 - \lambda_1}{\lambda_3 - \lambda_2} \frac{\lambda_4 - \lambda_2}{\lambda_4 - \lambda_1}.$$
 (2.22)

Moreover, the function

$$y(t) = -\frac{t}{1 + (1 - t)\hat{A}_4^{12}/\hat{A}_2^{12}} = \frac{1}{1 - \frac{1 - t}{t}\hat{A}_3^{12}/\hat{A}_2^{12}}$$
(2.23)

is the solution of the sixth Painlevé equation (1.3) with the parameters given by Eq.(1.7).

Proof. If the set $\{A_j\}$ is a solution of the system (2.9), (2.3), then the monodromy data of the function Ψ , which solves the corresponding Eq. (2.1), are independent of $\{\lambda_j\}$ and λ . Define the new variable

$$\mu = \frac{\lambda_3 - \lambda_1}{\lambda_3 - \lambda_2} \frac{\lambda - \lambda_2}{\lambda - \lambda_1} \tag{2.24}$$

and consider

$$\hat{\Psi} = G_1^{-1} \Psi C_1^{-1} \tag{2.25}$$

as a function of μ . In the complex μ -plane the function Φ has singularities only at the points 0, 1, t, and ∞ with the behavior prescribed by Eqs. (2.25) and (2.12): it is normalized at $\mu = \infty$ by the condition

$$\hat{\Psi} \underset{\mu \to \infty}{=} \left(I + \mathcal{O} \left(\mu^{-1} \right) \right) \mu^{\frac{1}{4}\sigma_3},$$

and its monodromy data are independent of $\{\lambda_j\}$. Such a function is uniquely defined and depends on $\{\lambda_j\}$ only via the cross-ratio t. It means that the logarithmic derivative,

$$\frac{d\hat{\Psi}}{d\mu}\hat{\Psi}^{-1} = \frac{\hat{A}_2}{\mu} + \frac{\hat{A}_3}{\mu - 1} + \frac{\hat{A}_4}{\mu - t} \stackrel{\text{def}}{=} \hat{A}(\mu), \tag{2.26}$$

and, in particular, the matrices

$$\hat{A}_j = \frac{t_j}{2} G_1^{-1} G_j \sigma_3 G_j^{-1} G_1$$

also depend on $\{\lambda_j\}$ only via the variable t. The matrices \hat{A}_j can be rewritten as follows,

$$\hat{A}_j = -\frac{t_j}{4}\hat{G}_j^{-1}\sigma_3\hat{G}_j,$$

where

$$\hat{G}_j = \begin{pmatrix} \det G_{j1}^{11} & \det G_{j1}^{12} \\ \det G_{j1}^{21} & \det G_{j1}^{22} \end{pmatrix}, \qquad \det \hat{G}_j = \det G_j \det G_1.$$
 (2.27)

To complete the proof one has to recall that according to [7] the function y(t), which solves the equation $\hat{A}^{12}(y) = 0$, where $A^{12}(\cdot)$ is the corresponding matrix element of $\hat{A}(\cdot)$ (see (2.26)), is the solution of the sixth Painlevé equation.

Remark 2.1 Proposition 2.4 is valid not only for the present case, when all coefficients t_j equal to $\frac{1}{2}$, but also in the case of arbitrary complex t_j . In the latter case the function y(t) (1.3) solves the sixth Painlevé equation with the coefficients:

$$\alpha = \frac{1}{2}(t_1 - 1)^2$$
, $\beta = -\frac{1}{2}t_2^2$, $\gamma = \frac{1}{2}t_3^2$, $\delta = \frac{1}{2}(1 - t_4^2)$.

The object playing the important role in applications of isomonodromy deformations in differential geometry and mathematical physics is the so-called tau function $\tau(\{\lambda_j\})$. We recall here the definition of the τ -function given in [7, 16, 8].

The Schlesinger equations (2.9), (2.10) can be rewritten in the Hamiltonian form,

$$\frac{dA_j}{d\lambda_k} = \{H_k, A_j\},\tag{2.28}$$

where the Poisson bracket is defined as follows,

$$\{(A_i)_{ab}, (A_j)_{cd}\} = \delta_{ij} ((A_i)_{ad} \delta_{cd} - (A_i)_{bc} \delta_{ad}), \qquad (2.29)$$

and the Hamiltonians are given by

$$H_j = \frac{1}{2} \underset{\lambda = \lambda_j}{\text{Res }} \operatorname{Tr} A^2(\lambda) = - \underset{\lambda = \lambda_j}{\text{Res }} \det A(\lambda) \equiv \sum_{i \neq j}^{2g+2} \frac{\operatorname{tr} A_j A_i}{\lambda_j - \lambda_i}.$$
 (2.30)

One proves that

$$\{H_k, H_j\} = 0, \qquad \frac{\partial H_k}{\partial \lambda_j} = \frac{\partial H_j}{\partial \lambda_k},$$
 (2.31)

which imply the compatibility of system (2.28). Taking into account the previous equations one can correctly define the τ -function $\tau \equiv \tau(\lambda_1, \ldots, \lambda_{2g+2})$ generating Hamiltonians H_i by

$$\frac{d}{d\lambda_i} \ln \tau = H_j, \tag{2.32}$$

which is holomorphic outside of the hyperplanes $\lambda_j = \lambda_i$, $i, j = 1, \dots, 2g + 2$.

3 Solutions of the Schlesinger System

Consider the hyperelliptic curve \mathcal{L} of genus g defined by the equation

$$w^{2} = \prod_{j=1}^{2g+2} (\lambda - \lambda_{j})$$
 (3.1)

with arbitrary non-coinciding $\lambda_j \in \mathbb{C}$ and the basic cycles (a_j, b_j) chosen according to figure 2.

Let us denote the fundamental polygon of \mathcal{L} by $\hat{\mathcal{L}}$. The basic holomorphic 1-forms on \mathcal{L} are given by

$$dU_k^0 = \frac{\lambda^{k-1}d\lambda}{w}, \qquad k = 1, \dots, g.$$
(3.2)

Let us define $g \times g$ matrices of a- and b-periods of these 1-forms by

$$\mathcal{A}_{kj} = \oint_{a_j} dU_k^0, \qquad \mathcal{B}_{kj} = \oint_{b_j} dU_k^0. \tag{3.3}$$

Then the holomorphic 1-forms

$$dU_k = \frac{1}{w} \sum_{j=1}^g (\mathcal{A}^{-1})_{kj} \lambda^{j-1} d\lambda \tag{3.4}$$

satisfy the normalization conditions $\oint_{a_i} dU_k = \delta_{jk}$.

The matrices \mathcal{A} and \mathcal{B} define the symmetric $g \times g$ matrix of b-periods of the curve \mathcal{L} :

$$\mathbf{B} = \mathcal{A}^{-1}\mathcal{B}$$
.

Let us now introduce the theta function with characteristic $[\mathbf{p}, \mathbf{q}]$ $(\mathbf{p} \in \mathbb{C}^g, \mathbf{q} \in \mathbb{C}^g)$ by the following series,

$$\Theta[\mathbf{p}, \mathbf{q}](\mathbf{z}|\mathbf{B}) = \sum_{\mathbf{m} \in \mathbb{Z}^g} \exp\{\pi i \langle \mathbf{B}(\mathbf{m} + \mathbf{p}), \mathbf{m} + \mathbf{p} \rangle + 2\pi i \langle \mathbf{z} + \mathbf{q}, \mathbf{m} + \mathbf{p} \rangle\},$$
(3.5)

for any $\mathbf{z} \in \mathbb{C}^g$. It possesses the following periodicity properties:

$$\Theta[\mathbf{p}, \mathbf{q}](\mathbf{z} + \mathbf{e}_j) = e^{2\pi i p_j} \Theta[\mathbf{p}, \mathbf{q}](\mathbf{z}), \tag{3.6}$$

$$\Theta[\mathbf{p}, \mathbf{q}](\mathbf{z} + \mathbf{B}\mathbf{e}_j) = e^{-2\pi i q_j} e^{-\pi i \mathbf{B}_{jj} - 2\pi i \mathbf{z}_j} \Theta[\mathbf{p}, \mathbf{q}](\mathbf{z}), \tag{3.7}$$

where

$$\mathbf{e}_j \equiv (0, \dots, 1, \dots, 0) \tag{3.8}$$

(1 stands in the jth place).

Denote the universal covering of \mathcal{L} by Γ . The multi-valued on \mathcal{L} , and single-valued on Γ , map $U(P) \in \mathbb{C}^g$ is defined by the contour integral $U_j(P) = \int_{\lambda_1}^P dU_j$. The vector of Riemann constants corresponding to our choice of the initial point of the map reads as follows [4]:

$$K = \frac{1}{2}\mathbf{B}(\mathbf{e}_1 + \dots + \mathbf{e}_g) + \frac{1}{2}(\mathbf{e}_1 + 2\mathbf{e}_2 \dots + g\mathbf{e}_g).$$
 (3.9)

The characteristic with components $\mathbf{p} \in \mathbb{C}^g/2\mathbb{C}^g$, $\mathbf{q} \in \mathbb{C}^g/2\mathbb{C}^g$ is called half-integer characteristic: the half-integer characteristics are in one-to-one correspondence with the half-periods $\mathbf{Bp} + \mathbf{q}$. If the scalar product $4\langle \mathbf{p}, \mathbf{q} \rangle$ is odd, then the related theta function is odd with respect to its argument \mathbf{z} and the characteristic $[\mathbf{p}, \mathbf{q}]$ is called odd, and if this scalar product is even, then the theta function $\Theta[\mathbf{p}, \mathbf{q}](\mathbf{z})$ is even with respect to \mathbf{z} and the characteristic $[\mathbf{p}, \mathbf{q}]$ is called even.

The odd characteristics which will be of importance for us in the sequel correspond to any given subset $S = \{\lambda_{i_1}, \dots, \lambda_{i_{g-1}}\}$ of g-1 arbitrary non-coinciding branch points. The odd half-period associated to the subset S is given by

$$\mathbf{Bp}^{S} + \mathbf{q}^{S} = U(\lambda_{i_1}) + \ldots + U(\lambda_{i_{q-1}}) - K. \tag{3.10}$$

Analogously, we shall be interested in the even half-periods which may be represented as follows,

$$\mathbf{B}\mathbf{p}^T + \mathbf{q}^T = U(\lambda_{i_1}) + \ldots + U(\lambda_{i_{g+1}}) - K, \tag{3.11}$$

where $T = \{\lambda_{i_1}, \dots, \lambda_{i_{q+1}}\}$ is an arbitrary subset of g+1 branch points.

Theorem 3.1 Let the 2×2 matrix-valued function $\Phi(P)$ be defined on the universal covering Γ of \mathcal{L} by the following formula,

$$\Phi(P) = \begin{pmatrix} \varphi(P) & \varphi(P^*) \\ \psi(P) & \psi(P^*) \end{pmatrix}, \tag{3.12}$$

where

$$\varphi(P) = \Theta[\mathbf{p}, \mathbf{q}](U(P) + U(P_{\varphi}))\Theta[\mathbf{p}^{S}, \mathbf{q}^{S}](U(P) - U(P_{\varphi})), \tag{3.13}$$

$$\psi(P) = \Theta[\mathbf{p}, \mathbf{q}](U(P) + U(P_{\psi}))\Theta[\mathbf{p}^S, \mathbf{q}^S](U(P) - U(P_{\psi})), \tag{3.14}$$

with arbitrary (possibly $\{\lambda_j\}$ -dependent) $P_{\varphi,\psi} \in \mathcal{L}$ and arbitrary constant characteristic $[\mathbf{p}, \mathbf{q}]$; * is the involution on \mathcal{L} interchanging the sheets. The odd theta characteristic $[\mathbf{p}^S, \mathbf{q}^S]$ corresponds to an arbitrary subset S of g-1 branch points via Eq. (3.10).

Then the function $\Phi(P)$ is holomorphic and invertible outside of the branch points $\lambda_1, \ldots, \lambda_{2g+2}$ and transforms as follows with respect to the tracing along the basic cycles of \mathcal{L} .

$$T_{a_j}[\Phi(P)] = \Phi(P)e^{2\pi i(p_j + p_j^S)\sigma_3}, \qquad T_{b_j}[\Phi(P)] = \Phi(P)e^{-2\pi i(q_j + q_j^S)\sigma_3}e^{-2\pi i\mathbf{B}_{jj} - 4\pi iU(P)},$$
(3.15)

where by T_l we denote the operator of analytic continuation along the contour l. Moreover, the function Φ has the following asymptotic expansion in the neighborhood of point λ_i :

$$\Phi(P) = \atop{\lambda \to \lambda_j} \left\{ F_j + O(\sqrt{\lambda - \lambda_j}) \right\} \begin{pmatrix} (\lambda - \lambda_j)^{1/2 + \delta_j} & 0 \\ 0 & (\lambda - \lambda_j)^{\delta_j} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix}, \quad (3.16)$$

with some λ -independent matrices F_j , $j = 1, \ldots, 2g + 2$; $\delta_j = 1$ for $\lambda_j \in S$ and $\delta_j = 0$ for $\lambda_j \notin S$.

Proof. Let us first check the announced monodromy properties of $\Phi(P)$ around the basic cycles of \mathcal{L} . From the periodicity properties of the theta function given by Eqs.(3.6), (3.7) we deduce the following transformation laws for φ :

$$T_{a_j}[\varphi(P)] = e^{2\pi i (p_j + p_j^S)} \varphi(P), \tag{3.17}$$

$$T_{b_j}[\varphi(P)] = e^{-2\pi i (q_j + q_j^S)} e^{-2\pi i \mathbf{B}_{jj} - 4\pi i U(P)} \varphi(P), \tag{3.18}$$

and the same transformation laws for ψ . Taking into account the action of the involution * on the basic cycles and holomorphic differentials,

$$a_j^* = -a_j, b_j^* = -b_j, dU_j(P^*) = -dU_j(P),$$
 (3.19)

we get the transformation laws for the function $\varphi(P^*)$,

$$T_{a_j}[\varphi(P^*)] = e^{-2\pi i(p_j + p_j^S)} \varphi(P^*),$$
 (3.20)

$$T_{b_j}[\varphi(P^*)] = e^{2\pi i (q_j + q_j^S)} e^{-2\pi i \mathbf{B}_{jj} - 4\pi i U(P)} \varphi(P^*), \tag{3.21}$$

which coincide with the transformation laws for the function $\psi(P^*)$. Altogether, this implies relations (3.15) for the function $\Phi(P)$.

The holomorphy of the function Φ follows from the holomorphy of the theta function. Let us show that det Φ does not vanish outside of the branch points λ_j . Since the transformations (3.15) preserve the positions of the zeros of det Φ , it makes sense to speak about the positions of the zeros of det Φ in the fundamental polygon $\hat{\mathcal{L}}$. First, notice that det $\Phi(P)$ vanishes at the branch points λ_j , where two columns of the matrix Φ coincide. Moreover, det Φ has at the points $\lambda_j \in S$ zeros of order 3. This can be seen if we rewrite the second theta function in Eq. (3.13) up to a non-vanishing exponential factor as

$$\Theta(U(P) - U(P_{\varphi}) - \sum_{S} U(\lambda_j) - K).$$

Thus we know altogether 3(g-1)+g+3=4g zeroes of $\det \Phi$ taking into account their multiplicities. To check that $\det \Phi$ does not vanish outside of λ_j , we integrate the function $\frac{\partial}{\partial \lambda} \ln \det \Phi(P)$ along the boundary of the fundamental polygon $\partial \hat{\mathcal{L}}$. From the transformation properties (3.15) we deduce

$$T_{a_j}[\det \Phi(P)] = \det \Phi(P), \qquad T_{b_j}[\det \Phi(P)] = e^{-4\pi i \mathbf{B}_{jj} - 8\pi i U_j(P)} \det \Phi(P).$$
 (3.22)

Now one can check that this integral equals 4g in the same way as in the standard calculation of the number of zeros of theta-function of dimension g [14]. Therefore det $\Phi(P)$ does not have any zeros outside of the branch points λ_i .

The form of the asymptotic expansion (3.16) is a direct consequence of the holomorphicity of φ and ψ , the structure (3.12) of the function Φ , and the previous discussion of the zeros of det Φ .

Starting from the function $\Phi(P)$ on Γ constructed in the Theorem 3.1, we shall now define a new function $\Psi(Q)$ on the universal covering X of $\mathbb{C} \setminus \{\lambda_1, \ldots, \lambda_{2g+2}\}$. Let us denote by $\Omega \subset \mathbb{C}$ an arbitrary neighborhood of ∞ on \mathbb{C} which does not overlap with the points λ_j and the projections of all basic cycles of \mathcal{L} on \mathbb{C} . Let us fix some sheet X_0 of X choosing the branch cuts between the points λ_j to lie outside of domain Ω . Let us also fix some sheet $\hat{\mathcal{L}}$ of the universal covering Γ of \mathcal{L} ; then $\hat{\mathcal{L}}$ will contain two non-intersecting copies of Ω . Choose one of them and denote by Ω_1 . The domain Ω_1 contains the point at infinity, which we call ∞^1 . Now we are in position to define

$$\Psi(\lambda \in \Omega) = \sqrt{\frac{\det \Phi(\infty^1)}{\det \Phi(\lambda)}} \Phi^{-1}(\infty^1) \Phi(\lambda)$$
 (3.23)

(by λ we denote the projection of $Q \in X$ as well as of $P \in \Gamma$ on \mathbb{C}). On the rest of X the function $\Psi(Q)$ is defined via the analytic continuation along the contours l_i (Fig.1).

Theorem 3.2 Let $\mathbf{p}, \mathbf{q} \in \mathbb{C}^g$ be an arbitrary set of 2g constants such that $[\mathbf{p}, \mathbf{q}]$ is not a half-integer characteristic. Then the function $\Psi(Q \in X)$ defined by (3.23), (3.12) is independent of the choice of the points $P_{\varphi,\psi} \in \mathcal{L}$ and the choice of the set $S = \{\lambda_{i_1}, \ldots, \lambda_{i_{g-1}}\}$. Moreover, Ψ is holomorphic outside of the branch points $\lambda_1, \ldots, \lambda_{2g+2}$, satisfies the normalization conditions $\det \Psi(\lambda) = 1$ and $\Psi(\lambda = \infty) = I$, and has the anti-diagonal monodromies M_j given by Eq. (2.17) along the contours l_j (Fig.1). The matrix elements of the monodromies (2.17) are given by the following expressions:

$$m_{1} = i, m_{2} = i(-1)^{g} \exp\{-2\pi i \sum_{k=1}^{g} p_{k}\},$$

$$m_{2j+1} = i(-1)^{g+1} \exp\{2\pi i q_{j} - 2\pi i \sum_{k=j}^{g} p_{k}\},$$

$$m_{2j+2} = i(-1)^{g} \exp\{2\pi i q_{j} - 2\pi i \sum_{k=j+1}^{g} p_{k}\},$$

$$(3.24)$$

for j = 1, ..., g, where p_j and q_j are components of the vectors \mathbf{p} and \mathbf{q} , respectively. The asymptotic expansion of $\Psi(Q)$ in the neighborhood of λ_j is of the form (2.12) with some G_j and

$$C_j = \frac{1}{\sqrt{2im_j}} \begin{pmatrix} 1 & im \\ -1 & im \end{pmatrix}. \tag{3.25}$$

Proof. The non-trivial part is to calculate the monodromies M_j of $\Psi(P)$ along the contours l_j .

Combining the transformations (3.15) of function Φ along the basic cycles of \mathcal{L} with the jumps of Φ ,

$$\Phi(P) \to \Phi(P)\sigma_1$$
,

on the branch cuts $[\lambda_{2j+1}, \lambda_{2j+2}]$, which follow directly from the definition (3.12), we come to the following relations:

$$\Psi(P)M_{2j+2}M_{2j+1} = \frac{T_{l_{2j+1}\circ l_{2j+2}}[\sqrt{\det\Phi(P)}]}{\sqrt{\det\Phi(P)}}\Psi(P)e^{2\pi i(p_j - p_j^S)\sigma_3},$$
 (3.26)

$$\Psi(P)M_{2j+1}M_{2j} = \frac{T_{l_{2j}\circ l_{2j+1}}[\sqrt{\det\Phi(P)}]}{\sqrt{\det\Phi(P)}}\Psi(P)e^{2\pi i(q_j-q_{j-1}+q_j^S-q_{j-1}^S)\sigma_3},$$
(3.27)

 $j = 1, \dots, g$. Furthermore, taking into account that

$$U(\lambda_1) = 0, \qquad U(\lambda_2) = \frac{1}{2} \sum_{k=1}^{g} \mathbf{e}_k,$$

$$U(\lambda_{2j+1}) = \frac{1}{2}\mathbf{B}\mathbf{e}_j + \frac{1}{2}\sum_{k=j}^{g}\mathbf{e}_k, \qquad U(\lambda_{2j+2}) = \frac{1}{2}\mathbf{B}\mathbf{e}_j + \frac{1}{2}\sum_{k=j+1}^{g}\mathbf{e}_k, \qquad j = 1, \dots, g, \quad (3.28)$$

we get

$$p_j^S = \frac{1}{2}(\delta_{2j+1} + \delta_{2j+2} + 1), \qquad q_{j+1}^S - q_j^S = \frac{1}{2}(\delta_{2j+2} + \delta_{2j+3} + 1),$$
 (3.29)

where δ_j are the same as in Eq. (3.16).

The function $\sqrt{\det \Phi(P)}$ transforms in the following way with respect to the tracing along the cycles l_i :

$$T_{l_{2j+1} \circ l_{2j+2}}[\sqrt{\det \Phi(P)}] = e^{\pi i(\delta_{2j+1} + \delta_{2j+2} + 1)} \sqrt{\det \Phi(P)}, \tag{3.30}$$

$$T_{l_{2j} \circ l_{2j+1}}[\sqrt{\det \Phi(P)}] = e^{\pi i (\delta_{2j+2} + \delta_{2j+3} + 1)} \sqrt{\det \Phi(P)}.$$
 (3.31)

To prove relations (3.30), (3.31) it is enough to notice that in the λ -plane the function $\sqrt{\det \Phi(P)}$ has at the point λ_j a zero of degree 3/4 if $\lambda_j \in S$ and zero of degree 1/4 if $\lambda_j \notin S$.

Altogether we get

$$M_{2j+2}M_{2j+1} = \exp\{2\pi i p_j \sigma_3\},$$

$$M_{2j+1}M_{2j} = \exp\{2\pi i (q_j - q_{j-1})\sigma_3\},$$

which imply (3.24) taking into account that $m_1 = i$ and the monodromy around infinity is trivial (2.18).

Now the independence of the function Ψ of the choice of the divisor S and the points $P_{\varphi,\psi}$ follows from the uniqueness of the solution to the Riemann-Hilbert problem with fixed monodromy data.

Existence of the local espansion (2.12) of the function $\Psi(Q)$ at the points λ_j follows from the related statement (3.16) for the function Φ which was proved in Theorem 3.1. The form (3.25) of the matrices C_j follows from the relation (2.11) between the matrices M_j and C_j .

Remark 3.1 The assumption made in Theorem 3.2 that $[\mathbf{p}, \mathbf{q}]$ does not coincide with any half-integer characteristic is nothing but the non-triviality condition, namely, if $[\mathbf{p}, \mathbf{q}]$ is a half-integer characteristic, all monodromies M_j become proportional to σ_1 : $M_j = \pm i\sigma_1$; therefore, they can be simultaneously diagonalized by the transformation

$$\Psi \rightarrow \tilde{\Psi} \equiv \Psi \left(\begin{array}{cc} 1 & 1 \\ 1 & -1 \end{array} \right),$$

The function $\tilde{\Psi}$ has diagonal monodromies $\pm i\sigma_3$, and, therefore, can be chosen to be diagonal itself. Thus, we are in the framework of the scalar Riemann-Hilbert problem: the related matrices A_j are diagonal, and, therefore, λ_j -independent, as follows from the Schlesinger equations.

By the special choice $P_{\varphi} = \infty^2$ and $P_{\psi} = \infty^1$ in the formulas of Theorem 3.1, we can simplify the previous expression for the function Ψ to get the following

Corollary 3.1 The function $\Psi(\lambda)$ defined by Eq. (3.23) may alternatively be represented as follows:

$$\Psi(\lambda \in \Omega) = \frac{1}{\sqrt{\det \Phi^{\infty}(\lambda)}} \Phi^{\infty}(\lambda), \tag{3.32}$$

where

$$\Phi^{\infty}(P) = \begin{pmatrix} \varphi^{\infty}(P) & \varphi^{\infty}(P^*) \\ \psi^{\infty}(P) & \psi^{\infty}(P^*) \end{pmatrix}, \tag{3.33}$$

$$\varphi^{\infty}(P) = \frac{\Theta[\mathbf{p}, \mathbf{q}](U(P) + U(\infty^2))\Theta[\mathbf{p}^S, \mathbf{q}^S](U(P) - U(\infty^2))}{\Theta[\mathbf{p}, \mathbf{q}](0)\Theta[\mathbf{p}^S, \mathbf{q}^S](-2U(\infty^2))},$$
(3.34)

$$\psi^{\infty}(P) = \frac{\Theta[\mathbf{p}, \mathbf{q}](U(P) + U(\infty^{1}))\Theta[\mathbf{p}^{S}, \mathbf{q}^{S}](U(P) - U(\infty^{1}))}{\Theta[\mathbf{p}, \mathbf{q}](0)\Theta[\mathbf{p}^{S}, \mathbf{q}^{S}](-2U(\infty^{1}))}.$$
 (3.35)

From the asymptotic expansions of the function $\Phi^{\infty}(P)$ at the points λ_j we can now construct solutions to the Schlesinger system.

Theorem 3.3 The solution to the Schlesinger system (2.9), (2.10) corresponding to the monodromy matrices (2.17), (3.24) is given by

$$A_{j} = \frac{1}{4} F_{j}^{\infty} \sigma_{3}(F_{j}^{\infty})^{-1}, \tag{3.36}$$

where

$$(F_j^{\infty})^{11} = \frac{\Theta[\mathbf{p}, \mathbf{q}](U(\lambda_j) + U(\infty^2))\Theta[\mathbf{p}^{S_j}, \mathbf{q}^{S_j}](U(\lambda_j) - U(\infty^2))}{\Theta[\mathbf{p}, \mathbf{q}](0)\Theta[\mathbf{p}^{S_j}, \mathbf{q}^{S_j}](-2U(\infty^2))},$$
(3.37)

$$(F_j^{\infty})^{12} = \sum_{l=1}^g \frac{\sum_{l=1}^g (\mathcal{A}^{-1})_{lk} \lambda_j^{l-1}}{\prod_{l \neq j} (\lambda_j - \lambda_l)^{1/2}}$$

$$\times \frac{\partial}{\partial z_k} \left\{ \frac{\Theta[\mathbf{p}, \mathbf{q}](\mathbf{z} + U(\infty^2))\Theta[\mathbf{p}^{S_j}, \mathbf{q}^{S_j}](\mathbf{z} - U(\infty^2))}{\Theta[\mathbf{p}, \mathbf{q}](0)\Theta[\mathbf{p}^{S_j}, \mathbf{q}^{S_j}](-2U(\infty^2))} \right\} (\mathbf{z} = U(\lambda_j)), \tag{3.38}$$

and $\partial/\partial z_k$ means the derivative of the theta function (3.5) with respect to its kth variable; matrix \mathcal{A} is given by Eq. (3.3); S_j are arbitrary 2g+2 sets of g-1 branch points λ_j satisfying the conditions $\lambda_j \notin S_j$. The solution (3.36) is independent of the choice of the sets S_j as long as these conditions are fulfilled.

The formulas for the matrix elements $(F_j^{\infty})^{21}$ and $(F_j^{\infty})^{22}$ may be obtained from the formulas for $(F_j^{\infty})^{11}$ and $(F_j^{\infty})^{12}$, respectively, by interchanging ∞^1 and ∞^2 .

Proof. In the neighborhood of the point λ_i we have

$$\varphi_j^{\infty}(P) = (F_j^{\infty})^{11} + \sqrt{\lambda - \lambda_j} (F_j^{\infty})^{12} + O(\lambda - \lambda_j), \tag{3.39}$$

$$\psi_j^{\infty}(P) = (F_j^{\infty})^{21} + \sqrt{\lambda - \lambda_j} (F_j^{\infty})^{22} + O(\lambda - \lambda_j), \tag{3.40}$$

with F_j given by Eqs. (3.37), (3.38); the functions $\varphi_j^{\infty}(P)$ and $\psi_j^{\infty}(P)$ are defined by Eqs. (3.13), (3.14), with $P_{\varphi} = \infty^2$, $P_{\psi} = \infty^1$, and $[\mathbf{p}^S, \mathbf{q}^S]$ substituted by $[\mathbf{p}^{S_j}, \mathbf{q}^{S_j}]$. Therefore,

$$\det \Phi_i^{\infty}(P) = \sqrt{\lambda - \lambda_j} \{ \det F_i^{\infty} + O(\lambda - \lambda_j) \}, \tag{3.41}$$

and

$$[\det \Phi_{j}^{\infty}(P)]^{-1/2}\varphi_{j}^{\infty}(P) = [\det F_{j}^{\infty}]^{-1}[(F_{j}^{\infty})^{11} + \sqrt{\lambda - \lambda_{j}}(F_{j}^{\infty})^{12} + O(\lambda - \lambda_{j})],$$

$$[\det \Phi_{j}^{\infty}(P)]^{-1/2}\psi_{j}^{\infty}(P) = [\det F_{j}^{\infty}]^{-1}[(F_{j}^{\infty})^{21} + \sqrt{\lambda - \lambda_{j}}(F_{j}^{\infty})^{22} + O(\lambda - \lambda_{j})].$$

We conclude that the matrices G_j , from the asymptotic expansions (2.12) of the function $\Psi(Q)$ at the points λ_j , are given by

$$G_j = (\det F_i^{\infty})^{-1} F_i^{\infty}, \tag{3.42}$$

which proves Eq. (3.36).

Remark 3.2 The matrices F_j^{∞} from Theorem 3.3 are related to the coefficients F_j of the asymptotic expansions (3.16) of function $\Phi(P)$ at the points λ_j as follows,

$$F_j^{\infty} = \Phi^{-1}(\infty^1) F_j,$$

Therefore, using Eq. (3.42), we get the following relation between the matrices F_j from the asymptotic expansions (3.16) of function $\Phi(P)$ and the matrices G_j from the asymptotic expansions (2.12) of function $\Psi(Q)$:

$$F_k^{-1}F_i\sigma_3F_i^{-1}F_k = G_k^{-1}G_i\sigma_3G_i^{-1}G_k, (3.43)$$

for any j and k.

4 Tau function for the Schlesinger System

Here we calculate the τ -function which corresponds to the solution (3.36), (3.37), (3.38) of the Schlesinger system. The remainder is devoted to the proof of the following main

Theorem 4.1 The τ -function corresponding to the solution (3.36), (3.37), (3.38) of the Schlesinger system (with arbitrary $\mathbf{p}, \mathbf{q} \in \mathbb{C}^g$ such that $[\mathbf{p}, \mathbf{q}]$ is not a half-integer characteristic) is given by

$$\tau = \Theta[\mathbf{p}, \mathbf{q}](0)(\det \mathcal{A})^{-1/2} \prod_{j \le k} (\lambda_j - \lambda_k)^{-1/8}, \tag{4.1}$$

where the $g \times g$ matrix A of a-periods of holomorphic 1-forms on \mathcal{L} is defined by Eq. (3.3).

Proof. According to the definition of the τ -function (2.32), (2.30), let us first calculate $\frac{1}{2} \text{tr}(\Psi_{\lambda} \Psi^{-1})^2$ for the function Ψ given by Eq. (3.23). We have

$$\frac{1}{2}\operatorname{tr}(\Psi_{\lambda}\Psi^{-1})^{2} \equiv -\det(\Psi_{\lambda}\Psi^{-1}) = -\frac{\det(\Phi_{\lambda})}{\det\Phi} + \frac{1}{4}\left(\frac{(\det\Phi)_{\lambda}}{\det\Phi}\right)^{2}.$$
 (4.2)

Together with the function Ψ , the function $\det(\Psi_{\lambda}\Psi^{-1})$ is independent of P_{φ} and P_{ψ} ; moreover, function Ψ does not undergo any modification if we multiply $\psi(P)$ with an arbitrary λ -independent factor C_{ψ} . So, we can choose the parameters P_{φ} , P_{ψ} , and C_{ψ} at our disposal to simplify the calculation. Our choice will be the following: first we put $C_{\psi} = \lambda_{\psi} - \lambda_{\varphi}$ (λ_{φ} denotes the projection of the point P_{φ} in the λ -plane) and then take the limit $P_{\psi} \to P_{\varphi}$. We get

$$\psi(P) = \varphi(P) + \frac{\partial \varphi(P)}{\partial \lambda_{\omega}}.$$
(4.3)

Since the function $\Psi(P)$ is independent of the remaining parameter P_{φ} , we can calculate $\det(\Psi_{\lambda}\Psi^{-1})$ assuming $P_{\varphi} = P$. Intermediate results of this calculation are as follows:

$$\frac{(\det \Phi)_{\lambda}}{\det \Phi} = 2 \frac{\partial}{\partial \lambda} \left\{ \ln \Theta[\mathbf{p}^S, \mathbf{q}^S](-2U(P)) \right\},\,$$

and

$$\begin{split} \frac{\det(\Phi)_{\lambda}}{\det\Phi} &= \frac{1}{\Theta[\mathbf{p},\mathbf{q}](0)} \frac{\partial^{2}}{\partial\lambda\partial\lambda_{\varphi}} \left\{ \Theta[\mathbf{p},\mathbf{q}](U(P) - U(P_{\varphi})) \right\}_{P_{\varphi} = P} \\ &+ \frac{1}{\Theta[\mathbf{p}^{S},\mathbf{q}^{S}](-2U(P))} \frac{\partial^{2}}{\partial\lambda\partial\lambda_{\varphi}} \left\{ \Theta[\mathbf{p}^{S},\mathbf{q}^{S}]((-U(P) - U(P_{\varphi})) \right\}_{P_{\varphi} = P}; \end{split}$$

therefore,

$$\frac{1}{2} \operatorname{tr}(\Psi_{\lambda} \Psi^{-1})^{2}(\lambda) = -\frac{\partial^{2}}{\partial \lambda \partial \lambda_{\varphi}} \left\{ \ln \Theta[\mathbf{p}^{S}, \mathbf{q}^{S}] ((-U(P) - U(P_{\varphi}))) \right\}_{P_{\varphi} = P}$$

$$-\frac{1}{\Theta[\mathbf{p}, \mathbf{q}](0)} \frac{\partial^{2}}{\partial \lambda \partial \lambda_{\varphi}} \left\{ \Theta[\mathbf{p}, \mathbf{q}] (-U(P) + U(P_{\varphi})) \right\}_{P_{\varphi} = P}.$$
(4.4)

To find the asymptotic expansion of this expression as $\lambda \to \lambda_j$ we shall use the well-known asymptotic expansion which is valid for any odd theta-characteristic [$\mathbf{p}^S, \mathbf{q}^S$]:

$$\frac{\partial^2}{\partial x(P_1)\partial x(P_2)} \left\{ \ln \Theta[\mathbf{p}^S, \mathbf{q}^S](U(P_1) - U(P_2)) \right\} = \frac{1}{(x(P_1) - x(P_2))^2} + F(P) + o(1) \quad (4.5)$$

as $P_1, P_2 \to P$, where x is a local parameter in the neighborhood of P. The function F(P) is independent of the choice of the set S; it is given by the following expression ([4], p.20),,

$$F(P) = \frac{1}{6} \{\lambda, x\}(P) + \frac{1}{16} \left(\frac{d}{dx} \ln \prod_{k=1}^{g+1} \frac{\lambda - \lambda_{i_k}}{\lambda - \lambda_{j_k}}\right)^2 (P)$$
$$- \sum_{i,j=1}^g \frac{\partial^2}{\partial z_i \partial z_j} \Theta[\mathbf{p}^T, \mathbf{q}^T](0) \frac{dU_i}{dx} (P) \frac{dU_j}{dx} (P), \tag{4.6}$$

where $\{\lambda, x\}$ denotes the Schwarzian derivative of λ with respect to x,

$$\frac{\lambda'''}{\lambda'} - \frac{3}{2} \left(\frac{\lambda''}{\lambda'}\right)^2 ,$$

and $[\mathbf{p}^T, \mathbf{q}^T]$ is an even characteristic corresponding to an arbitrary set $T \equiv \{\lambda_{i_1}, \dots, \lambda_{i_{g+1}}\}$ of g+1 branch points via Eq. (3.11). The remaining g+1 branch points are denoted by $\lambda_{j_1}, \dots, \lambda_{j_{g+1}}$. Expression (4.6) is independent of the choice of the set T.

Applying Eq. (4.6) for $P = \lambda_j$ we get the following asymptotic expansion,

$$\frac{1}{2}\operatorname{tr}(\Psi_{\lambda}\Psi^{-1})^{2}(\lambda) \underset{\lambda \to \lambda_{j}}{=} \frac{1}{16(\lambda - \lambda_{j})^{2}} + \frac{H_{j}}{\lambda - \lambda_{j}} + O(1), \tag{4.7}$$

where

$$H_{j} = \frac{1}{8} \sum_{k \neq j} \frac{n_{j} n_{k}}{\lambda_{j} - \lambda_{k}} - \frac{1}{4\Theta[\mathbf{p}^{T}, \mathbf{q}^{T}](0)} \sum_{l,k=1}^{g} \frac{\partial^{2} \Theta[\mathbf{p}^{T}, \mathbf{q}^{T}]}{\partial z_{l} \partial z_{k}} (0) \frac{dU_{l}}{dx_{j}} (\lambda_{j}) \frac{dU_{k}}{dx_{j}} (\lambda_{j})$$

$$+ \frac{1}{4\Theta[\mathbf{p}, \mathbf{q}](0)} \sum_{l,k=1}^{g} \frac{\partial^{2} \Theta[\mathbf{p}, \mathbf{q}]}{\partial z_{l} \partial z_{k}} (0) \frac{dU_{l}}{dx_{j}} (\lambda_{j}) \frac{dU_{k}}{dx_{j}} (\lambda_{j}), \tag{4.8}$$

and $x_j \equiv \sqrt{\lambda - \lambda_j}$; $n_k = 1$ for $\lambda_k \in T$ and $n_k = -1$ for $\lambda_k \notin T$. Now, to integrate Eqs. (2.32), we have to use the heat equations

$$\frac{\partial^2 \Theta[\mathbf{p}, \mathbf{q}](\mathbf{z}|\mathbf{B})}{\partial z_l \partial z_k} = 4\pi i \frac{\partial \Theta[\mathbf{p}, \mathbf{q}](\mathbf{z}|\mathbf{B})}{\partial \mathbf{B}_{lk}}$$
(4.9)

valid for theta functions with arbitrary characteristic [p, q], and the following

Lemma 4.1 The dependence of the matrix of b-periods on the branch points is described by the following equations,

$$\frac{\partial \mathbf{B}_{kl}}{\partial \lambda_j} = \pi i \frac{dU_l}{dx_j} (\lambda_j) \frac{dU_k}{dx_j} (\lambda_j). \tag{4.10}$$

Proof. The dependence of the non-normalized 1-forms dU_k^0 (3.2) on λ_j is

$$\frac{\partial}{\partial \lambda_i} \{ dU_k^0(\lambda) \} = \frac{1}{2(\lambda - \lambda_i)} dU_k^0(\lambda).$$

Now, calculation of the integral

$$\oint_{\partial \hat{\mathcal{L}}} U_l^0(\lambda) \frac{\partial}{\partial \lambda_j} dU_k^0(\lambda) = \oint_{\partial \hat{\mathcal{L}}} \frac{1}{2(\lambda - \lambda_j)} U_l^0(\lambda) dU_k^0(\lambda)$$

by means of the residue theorem gives the following result:

$$\pi i \frac{dU_l^0}{dx_j}(\lambda_j) \frac{dU_k^0}{dx_j}(\lambda_j) \equiv \pi i \left[\mathcal{A} \frac{dU_l}{dx_j}(\lambda_j) \frac{dU_k}{dx_j}(\lambda_j) \mathcal{A}^t \right]_{kl}.$$

On the other hand, standard arguments used, for example, in the proof of the Riemann bilinear identities [6], show that the same integral equals

$$\sum_{m=1}^{g} \mathcal{A}_{lm} \frac{\partial \mathcal{B}_{km}}{\partial \lambda_j} - \frac{\partial \mathcal{A}_{km}}{\partial \lambda_j} \mathcal{B}_{lm};$$

therefore.

$$\frac{\partial \mathcal{B}}{\partial \lambda_j} \mathcal{A}^t - \frac{\partial \mathcal{A}}{\partial \lambda_j} \mathcal{B}^t = \pi i \mathcal{A} \frac{dU_l}{dx_j} (\lambda_j) \frac{dU_k}{dx_j} (\lambda_j) \mathcal{A}^t,$$

which leads to the statement of the lemma (4.10) after taking into account the symmetry of the matrix $\mathbf{B} \equiv \mathcal{A}^{-1}\mathcal{B}$.

Now, using Eqs. (4.8), (4.9), and (4.10), we can rewrite the Hamiltonians H_j as follows:

$$H_j \equiv \frac{\partial}{\partial \lambda_j} \ln \tau = \frac{1}{8} \sum_{k \neq j} \frac{n_j n_k}{\lambda_j - \lambda_k} + \frac{\partial}{\partial \lambda_j} \ln \left\{ \frac{\Theta[\mathbf{p}, \mathbf{q}](0)}{\Theta[\mathbf{p}^T, \mathbf{q}^T](0)} \right\}.$$

Finally, applying the classical Thomae formula [18, 14]

$$\Theta^{4}[\mathbf{p}^{T}, \mathbf{q}^{T}](0) = \pm (\det \mathcal{A})^{2} \prod_{l < k, \ l, k=1}^{g+1} (\lambda_{i_{l}} - \lambda_{i_{k}}) \prod_{l < k, \ l, k=1}^{g+1} (\lambda_{j_{l}} - \lambda_{j_{k}}),$$

we get the τ -function in the form (4.1) up to multiplication by an arbitrary $\{\lambda_j\}$ independent constant of integration. The ambiguity in the choice of this constant allows,
in particular, to arbitrarily choose the branch cuts in the formula (4.1).

5 Elliptic Case and Painlevé VI Equation

In this section we are going to show how the solution of the Painlevé VI equation in terms of elliptic functions can be derived from the results of the previous sections.

Put g=1. Then the equation of the curve \mathcal{L} is given by

$$w^{2} = (\lambda - \lambda_{1})(\lambda - \lambda_{2})(\lambda - \lambda_{3})(\lambda - \lambda_{4}). \tag{5.1}$$

The matrix of b-periods, **B**, turns into the module σ and $\Theta[\mathbf{p}^S, \mathbf{q}^S]$ becomes the Jacobi theta-function ϑ_1 ; to shorten all the formulas we shall denote $\Theta[\mathbf{p}, \mathbf{q}]$ by $\vartheta_{p,q}$.

Parameters m_j of the monodromy matrices are, according to (3.24), given by

$$m_1 = i$$
, $m_2 = -ie^{-2\pi ip}$, $m_3 = ie^{2\pi i(q-p)}$, $m_4 = -ie^{2\pi iq}$.

The formulas (3.13) and (3.14) now read as follows

$$\varphi(P) = \vartheta_{p,q}(U(P) + u_{\varphi})\vartheta_1(U(P) - u_{\varphi}), \tag{5.2}$$

$$\psi(P) = c_{\psi}\vartheta_{p,q}(U(P) + u_{\psi})\vartheta_1(U(P) - u_{\psi}), \tag{5.3}$$

where $u_{\varphi,\psi} \equiv U(P_{\varphi,\psi}) \in \mathbb{C}$ are arbitrary parameters, and, in analogy to the previous section, we introduced an arbitrary multiplier $c_{\psi}(\{\lambda_j\})$ which obviously does not influence the function $\Psi(\lambda)$.

Again, since the function $\Psi(\lambda)$ does not depend on c_{ψ} , u_{φ} and u_{ψ} , we can freely fix these parameters to simplify our calculations. First, it is convenient to put $u_{\varphi} = 0$ (i.e., $P_{\varphi} = \lambda_1$), which leads to

$$\varphi(P) = \vartheta_{p,q}(U(P))\vartheta_1(U(P)). \tag{5.4}$$

The most convenient choice for the parameters of the function ψ is the following: we put $c_{\psi} = u_{\psi}^{-1}$ and take the limit $u_{\psi} \to 0$. Then we get

$$\psi(P) = \varphi(P) + \frac{\partial \varphi(P)}{\partial u_{\varphi}} (u_{\varphi} = 0), \tag{5.5}$$

and the components of matrices F_j from Eq. (3.16) are given by

$$F_{j}^{11} = \vartheta_{p,q}(u_{j})\vartheta_{1}(u_{j}),$$

$$F_{j}^{12} = f_{j}\{\vartheta'_{p,q}(u_{j})\vartheta_{1}(u_{j}) + \vartheta_{p,q}(u_{j})\vartheta'_{1}(u_{j})\},$$

$$F_{j}^{21} = F_{j}^{11} + \vartheta'_{p,q}(u_{j})\vartheta_{1}(u_{j}) - \vartheta_{p,q}(u_{j})\vartheta'_{1}(u_{j}),$$

$$F_{j}^{22} = F_{j}^{12} + f_{j}\{\vartheta''_{p,q}(u_{j})\vartheta_{1}(u_{j}) - \vartheta_{p,q}(u_{j})\vartheta''_{1}(u_{j})\}.$$
(5.6)

(5.6)

Here

$$f_j \equiv \left\{ \prod_{l \neq j} (\lambda_j - \lambda_l)^{1/2} \oint_a \frac{d\lambda}{\sqrt{(\lambda - \lambda_1) \dots (\lambda - \lambda_4)}} \right\}^{-1}, \tag{5.7}$$

and

$$u_1 = 0, \quad u_2 = \frac{1}{2}, \quad u_3 = \frac{1}{2} + \frac{\sigma}{2}, \quad u_4 = \frac{\sigma}{2}.$$
 (5.8)

In particular, for j = 1 we have

$$F_1^{11} = 0, F_1^{21} = \vartheta_{p,q}(0)\vartheta_1'(0), F_1^{12} = F_1^{22} = f_1F_1^{21}. (5.9)$$

In accordance with Eqs. (3.43), (2.23), to obtain the solution of the sixth Painlevé equation we have to calculate the (12) elements of the matrices

$$\hat{A}_j = \frac{1}{4} F_1^{-1} F_j \sigma_3 F_j^{-1} F_1, \qquad j = 2, 3, 4$$
(5.10)

(obviously $\hat{A}_1 = I$). Substitution of the matrix elements (5.6) into Eq. (5.10) leads to the following result:

$$\hat{A}_{j}^{12} = -f_{1} \frac{((\ln \vartheta_{p,q})' - (\ln \vartheta_{1})')(\vartheta_{p,q}'' \vartheta_{p,q} - \vartheta_{1}'' / \vartheta_{1})}{(\ln \vartheta_{p,q})'' - (\ln \vartheta_{1})''} (z = u_{j}), \tag{5.11}$$

where ϑ' denotes for $\partial \vartheta(z|\sigma)/\partial z$. Finally, choosing $\lambda_1 = \infty$, $\lambda_2 = 0$, $\lambda_3 = 1$, and $\lambda_4 = t$, and making use of the "heat" equation for the theta-function,

$$\frac{\partial \vartheta_{p,q}(z,\sigma)}{\partial \sigma} = \frac{1}{4\pi i} \frac{\partial^2 \vartheta_{p,q}(z,\sigma)}{\partial z^2},$$

we get, according to Eq. (2.23), the following

Theorem 5.1 The function

$$y = -\frac{t}{1 + (1 - t)y_1} \quad , \tag{5.12}$$

where t is the cross-ratio of the points $\{\lambda_i\}$ given by Eq. (2.22), and

$$y_{1} = \frac{\frac{\partial}{\partial z} \ln \frac{\partial}{\partial z} \ln \{\vartheta_{p,q}/\vartheta_{1}\}(\frac{1}{2}) \frac{\partial}{\partial \sigma} \ln \{\vartheta_{p,q}/\vartheta_{1}\}(\frac{\sigma}{2})}{\frac{\partial}{\partial z} \ln \frac{\partial}{\partial z} \ln \{\vartheta_{p,q}/\vartheta_{1}\}(\frac{\sigma}{2}) \frac{\partial}{\partial \sigma} \ln \{\vartheta_{p,q}/\vartheta_{1}\}(\frac{1}{2})} .$$
 (5.13)

where $p,q \in \mathbb{C}$ are arbitrary constants such that $[p,q] \neq [1/2,0]$ and $[p,q] \neq [0,1/2]$, solves the sixth Painlevé equation (1.3), with coefficients (1.7). Here the module σ of elliptic curve \mathcal{L} is chosen such that $t = \theta_4^4(0|\sigma)/\theta_2^4(0|\sigma)$.

Expression (5.13) is a combination of derivatives of the function $\ln \frac{\vartheta_{p,q}}{\vartheta_1}$ with respect to both arguments of the theta functions.

One more representation for solution (5.12) of sixth Painlevé equation may be obtained by using the following relation between y(t) and the τ -function, $\tau(t)$, valid for $t_i = \frac{1}{2}$:

$$y(t) = t - t(t-1) \left[D\left(\frac{\frac{d}{dt}D(\tau)}{\frac{d}{dt}D\left(\sqrt[8]{t(t-1)\tau}\right)}\right) + \frac{t(t-1)}{D^2\left(\sqrt[8]{t(t-1)\tau}\right)} \right]^{-1}$$
 (5.14)

where operator D acts on functions f(t) as follows: $D(f) \equiv \frac{d}{dt} \ln f$. The τ -function for the g = 1 case can be obtained from the general formula (4.1) simply by assuming that $\lambda_1, \ldots, \lambda_4$ coincide with 0, 1, t, and ∞ , respectively. Then up to an arbitrary constant we get

$$\tau(t) = \frac{\theta_{p,q}(0|\sigma)}{\sqrt[8]{t(t-1)}} \left[\int_0^1 \frac{d\lambda}{\sqrt{\lambda(\lambda-1)(\lambda-t)}} \right]^{-\frac{1}{2}}.$$

Remark 5.1 It seems that it is not easy to check directly (by applying appropriate identities for the theta functions) the coincidence of the different forms of the same solution (5.13), (5.14). It is also not easy to check directly coincidence of our formulas to other forms of this solution given by Okamoto (A.6) and Hitchin (A.7). However, we can explicitly see the relationship of our construction to the construction by Hitchin on the level of the functions φ and ψ from Theorem 3.1, namely, the choice of the rows of the function Φ made in [5] corresponds to the choice $u_{\varphi} \equiv -\frac{1}{2}(p\sigma + q) + \frac{\sigma+1}{4}$. The variable c from [5] is given by $-u_{\varphi}w_1$, where w_1 is the first full elliptic integral on \mathcal{L} . The parameter u_{ψ} is fixed in [5] to coincide with one of the zeros of the Weierstrass \wp -function, $\wp[w_1(U(P) + u_{\varphi})]$, with the periods w_1 and $w_2 = w_1\sigma$. Constants c_1 and c_2 from [5] are related to our p and q as follows: $c_1 = p + \frac{1}{2}$, $c_2 = q + \frac{1}{2}$.

Remark 5.2 Here we discussed only generic two-parametric family of elliptic solutions of Painlevé 6 equation with coefficients (1.7), which corresponds to monodromy matrices (2.19). Additional one-parametric family of solutions corresponding to monodromy matrices (2.20) was constructed in [5].

A Elliptic Solutions of the Sixth Painlevé Equation

In his studies of the Painlevé equations K. Okamoto has shown [15] that the function y = y(t), the general solution of the sixth Painlevé equation, (1.3), can be explicitly written in terms of the elliptic functions provided the set of the parameters satisfies one of the following conditions:

$$t_i \in \mathbb{Z}, \quad t_1 + \ldots + t_4 \in 2\mathbb{Z},$$
 (A.1)

or

$$t_i + \frac{1}{2} \in \mathbb{Z}. \tag{A.2}$$

The major ingredients of the Okamoto's construction are:

1) The so-called Picard solution,

$$y_0(t) = \tilde{\wp}(c_1\omega_1(t) + c_2\omega_2(t)), \tag{A.3}$$

of Eq. (1.3) with the coefficients:

$$\alpha = 0, \quad \beta = 0, \quad \gamma = 0, \quad \delta = \frac{1}{2}.$$
 (A.4)

In Eq. (A.3) $\tilde{\wp}(\cdot)$ is the elliptic function satisfying the equation, $\tilde{\wp}'^2 = 4\tilde{\wp}(\tilde{\wp} - 1)(\tilde{\wp} - t)$, with the primitive periods $2\omega_1(t)$ and $2\omega_2(t)$; $c_1, c_2 \in \mathbb{C}$ are the constants of integration, so that the function y(t) is the general solution.

2) The subgroup of transformations of solutions of Eq. (1.3) which acts on the space of coefficients $\{t_j\}$ as: a)reflections: for any $j=1,\ldots,4$ there is a transformation which transforms $t_j \to -t_j$ and leaves all t_k for $k \neq j$ unchanged; b) permutations of the set $\{t_j\}$; c) the shifts: $t_j \mapsto t_j + n_j$, where $\sum_{j=1}^4 n_j = 0 \pmod{2}$.

3) More nontrivial transformation,

$$\mathbf{O}: (t_1, t_2, t_3, t_4) \leftrightarrow \left(\frac{t_1 + t_2 - t_3 - t_4}{2}, \frac{t_1 + t_2 + t_3 + t_4}{2}, \frac{-t_1 + t_2 + t_3 - t_4}{2}, \frac{-t_1 + t_2 - t_3 + t_4}{2}\right). (A.5)$$

It is important to mention that all the transformations described above, as well as their inversions, are given by explicit formulas, so that "new" solutions can be explicitly written in terms of the "old" ones as rational functions of the "old" solution and its derivative (see [15]). In particular, the solution of Eq. (1.3) with the coefficients (1.7) obtained via the Okamoto's transformations reads

$$y(t) = y_0 + \frac{y_0^2(y_0 - 1)(y_0 - t)}{t(t - 1)y_0' - y_0(y_0 - 1)},$$
(A.6)

where $y_0 = y_0(t)$ is given by Eq. (A.3).

N. Hitchin, in the work [5] devoted to the study of SU(2)-invariant anti-self-dual Einstein metrics, rediscovered the case (1.7) of integrability of Eq. (1.3) in elliptic functions. He got the following representation for the solution (A.6) in the parametric form,

$$y_{1}(\sigma) = \frac{\theta_{1}^{"''(0)}}{3\pi^{2}\theta_{4}^{4}(0)\theta_{1}^{\prime}(0)} + \frac{1}{3}\left(1 + \frac{\theta_{3}^{4}(0)}{\theta_{4}^{4}(0)}\right) + \frac{\theta_{1}^{"''(\nu)\theta_{1}(\nu) - 2\theta_{1}^{"}(\nu)\theta_{1}^{\prime}(\nu) + 2\pi i c_{1}(\theta_{1}^{"}(\nu)\theta_{1}(\nu) - \theta_{1}^{\prime^{2}(\nu))}}{2\pi^{2}\theta_{4}^{4}(0)\theta_{1}(\nu)(\theta_{1}^{\prime}(\nu) + \pi i c_{1}\theta_{1}(\nu))}, (A.7)$$

$$t(\sigma) = \frac{\theta_{3}^{4}(0)}{\theta_{4}^{4}(0)}, \qquad \nu = c_{1}\sigma + c_{2},$$

where $\theta_k(\cdot) = \theta(\cdot|\sigma), \ k = 1, \dots, 4$, are the Jacobi theta functions [19].

Yu. I. Manin [13] noticed that the well-known uniformization of the Eq. (1.3) in terms of the Weierstrass \wp -function can be further converted into the beautiful form:

$$y(\sigma) = \frac{\wp(z(\sigma), \sigma) - e_1(\sigma)}{e_2(\sigma) - e_1(\sigma)}, \qquad t(\sigma) = \frac{e_3(\sigma) - e_1(\sigma)}{e_2(\sigma) - e_1(\sigma)},$$

$$e_j(\sigma) = \wp(\frac{1}{2}T_j, \sigma), \qquad (T_1, T_2, T_3, T_4) \equiv (0, 1, \sigma, 1 + \sigma),$$

$$\frac{d^2z}{d\sigma^2} = \frac{1}{(2\pi i)^2} \sum_{j=1}^4 \alpha_j \wp'(z + \frac{T_j}{2}, \sigma), \quad (\alpha_1, \alpha_2, \alpha_3, \alpha_4) \equiv (\alpha, -\beta, \gamma, \frac{1}{2} - \delta), \quad (A.8)$$

where $\wp(\cdot,\sigma)$ is the Weierstrass elliptic function with the primitive periods 2 and 2σ ; $\wp'(\cdot,\sigma)$ denotes partial derivative of \wp -function with respect to its first argument. By applying to Eq. (A.8) the Landin transform for the Weierstrass elliptic functions Manin found a new transformation for solutions of Eq. (1.3). In terms of the Manin variables, z and σ this transformation reads: if $z(\sigma)$ is any solution of Eq. (A.8) with the coefficients $\alpha_1 = \alpha_3$, $\alpha_2 = \alpha_4$, then $z(2\sigma)$ is the solution of Eq. (A.8) for $\alpha_1^{new} = 4\alpha_1$, $\alpha_2^{new} = 4\alpha_2$, $\alpha_3^{new} = \alpha_4^{new} = 0$. The converse statement is, of course, also true. Schematically, for the constants, t_i (1.4), we can write:

$$\mathbf{M}: (t_1, t_2, t_3 = t_1 - 1, t_4 = t_2) \leftrightarrow (2t_1 - 1, 2t_2, 0, 0).$$
 (A.9)

In the case (A.4) the Manin form of the sixth Painlevé equation (A.8) immediately reproduces the Picard solution (A.3). In terms of the parameters t_j Eqs. (A.4) read, $t_1 = 1$, $t_2 = 0$, $t_3 = 0$, and $t_4 = 0$. After the permutation we get the set $t_1 = 0$, $t_2 = 1$, $t_3 = 0$, and $t_4 = 0$, therefore, by setting the formal monodromies $t_1 = \frac{1}{2}$, $t_2 = -\frac{1}{2}$ in the r.-h. s. of (A.9) and choosing the left arrow in Manin transformation \mathbf{M} , one finds the second basic case of the integrability (1.7). The corresponding explicit formula can be written as the composition of the transformation corresponding to the permutation [15] and \mathbf{M} .

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